

QUASI-COMMUTING FAMILIES OF QUANTUM MINORS

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In [7] a combinatorial criterion for quasi-commutativity is established for pairs of quantum Plücker coordinates in the quantized coordinate algebra $\mathbb{C}_q[\mathcal{F}]$ of the flag variety of type A . This paper attempts to generalize these results by producing necessary and sufficient conditions for pairs of quantum minors in the quantized coordinate algebra $\mathbb{C}_q[Mat_{k \times m}]$ to quasi-commute. In addition we study the combinatorics of maximal (by inclusion) families of pairwise quasi-commuting quantum minors and pose relevant conjectures.

1. INTRODUCTION

Let $\mathbb{C}_q[Mat_{k \times m}]$ be the q -deformation of the coordinate ring of the space of $k \times m$ complex matrices where $k \leq m$. This is the $\mathbb{C}(q)$ -algebra with unity generated by indeterminates $x_{i,j}$ for $i \in [1 \dots k]$ and $j \in [1 \dots m]$ subject to the Faddeev-Reshetikhin-Takhtadzhyan relations [2]:

$$\begin{aligned}
 x_{s,t}x_{i,j} &= q \, x_{i,j}x_{s,t} && \text{if either } s > i \text{ and } t = j \\
 &&& \text{or } s = i \text{ and } t > j \\
 x_{s,t}x_{i,j} &= x_{i,j}x_{s,t} && \text{if } s > i \text{ and } t < j \\
 x_{s,t}x_{i,j} &= x_{i,j}x_{s,t} + (q - q^{-1}) x_{i,t}x_{s,j} && \text{if } s > i \text{ and } t > j
 \end{aligned}$$

In this paper we shall be concerned with a special family of elements $\Delta_{I,J} \in \mathbb{C}_q[Mat_{k \times m}]$ indexed by pairs of non-empty subsets I and J of $[1 \dots k]$ and $[1 \dots m]$ respectively with $|I| = |J| = l$. They are defined by:

$$\Delta_{I,J} := \sum_{\sigma \in S_l} (-q)^{-l(\sigma)} x_{i_1, j_{\sigma(1)}} \cdots x_{i_l, j_{\sigma(l)}} ,$$

where $I = \{i_1 < \cdots < i_l\}$, $J = \{j_1 < \cdots < j_l\}$, and $l(\sigma)$ is the length of the l -permutation σ . The element $\Delta_{I,J}$ is the q -deformation of the classical determinant and for this reason we call the $\Delta_{I,J}$'s quantum minors.

Definition 1. Two quantum minors $\Delta_{A,B}$ and $\Delta_{C,D}$ quasi-commute if $\Delta_{C,D}\Delta_{A,B} = q^c \Delta_{A,B}\Delta_{C,D}$ for some integer c . The integer c is uniquely determined by $\Delta_{A,B}$ and $\Delta_{C,D}$ and we will denote its value by the symbol $c(\Delta_{A,B} \mid \Delta_{C,D})$. Note that $c(\Delta_{C,D} \mid \Delta_{A,B}) = -c(\Delta_{A,B} \mid \Delta_{C,D})$ for any quasi-commuting pair.

We can now state the central problems we will address in this paper, namely:

Problem 1. Find necessary and sufficient conditions for two quantum minors $\Delta_{A,B}$ and $\Delta_{C,D}$ to quasi-commute. In addition, explicitly compute $c(\Delta_{A,B} \mid \Delta_{C,D})$ in terms of A , B , C , and D .

Problem 2. Find a combinatorial mechanism which will describe and produce all maximal (by inclusion) families of pairwise quasi-commuting quantum minors.

Problems 1 and 2 are motivated by the study of dual canonical bases for quantum groups of type A . It is conjectured in [1], and partially proved in [8], that products of quasi-commuting quantum minors constitute a part of the dual canonical basis for the quantum group $\mathbb{C}_q[GL(n, \mathbb{C})]$. Problem 2 is also motivated by the study of total positivity as described in [3] and [4].

Problem 1 is resolved using techniques developed in [7]. Ostensibly Problem 1 is more general than its counterpart in [7] which only addresses the quantum flag variety. Nevertheless we demonstrate in this paper that Problem 1 can be reduced to a special case of the problem treated in [7] - namely the problem of determining when two quantum Plücker coordinates of the corresponding quantum Grassmannian quasi-commute. The criterion for quasi-commutativity is described in terms of the notion of "weak separability" as put forth in [7].

Definition 2. Given two subsets I and J of $[1 \dots n]$ we write $I \prec J$ if $i < j$ for all $i \in I$ and all $j \in J$. We say I and J are weakly separated if at least one of the following two conditions holds:

1. $|I| \geq |J|$ and $J - I$ can be partitioned into a disjoint union $J - I = J' \sqcup J''$ so that $J' \prec I - J \prec J''$.
2. $|J| \geq |I|$ and $I - J$ can be partitioned into a disjoint union $I - J = I' \sqcup I''$ so that $I' \prec J - I \prec I''$.

We associate to any pair of subsets $A \subset [1 \dots k]$ and $B \subset [1 \dots m]$ of equal size the subset $S(A, B) \subset [1 \dots k + m]$ of size k defined as follows:

$$S(A, B) = \left\{ b + k \mid b \in B \right\} \sqcup [1 \dots k] - w_0(A)$$

where w_0 is the order reversing permutation of $[1 \dots k]$. Problem 1 is settled by the following two Theorems:

Theorem 1. The quantum minors $\Delta_{A,B}$ and $\Delta_{C,D}$ in $\mathbb{C}_q[Mat_{k,m}]$ quasi-commute if and only if $S(A, B)$ and $S(C, D)$ are weakly separated subsets of $[1 \dots m + k]$.

Theorem 2. *Suppose $I = S(A, B)$ and $J = S(C, D)$ are weakly separated subsets of $[1 \dots m + k]$ satisfying case 1 in Definition 2. Then*

$$c(\Delta_{A,B} \mid \Delta_{C,D}) = |J''| - |J'| + |A| - |C|.$$

In proving Theorems 1 and 2 we use a quantum analogue of the well known embedding of $Mat_{k \times m}$ as an affine chart in the Grassmannian $\mathbb{G}_{k,k+m}$; this embedding sends a $k \times m$ matrix $(x_{i,j})$ to the row space of the $k \times (k + m)$ matrix

$$\begin{pmatrix} 0 & & & 1 & x_{1,1} & \cdot & \cdot & \cdot & \cdot & x_{1,m} \\ & & & -1 & \cdot & & & & & \cdot \\ & & \cdot & & \cdot & & & & & \cdot \\ & & & & \cdot & & & & & \cdot \\ & & \cdot & & \cdot & & & & & \cdot \\ & & & & \cdot & & & & & \cdot \\ (-1)^{k-1} & & & 0 & x_{k,1} & \cdot & \cdot & \cdot & \cdot & x_{k,m} \end{pmatrix}$$

The corresponding quantum analogue is an embedding of $\mathbb{C}_q[Mat_{k \times m}]$ into the quantized coordinate ring $\mathbb{C}_q[\mathbb{G}_{k,k+m}]$ - the so called quantum Grassmannian as defined in [10]. This embedding allows us to reduce questions about quantum minors to corresponding questions about quantum Plücker coordinates.

Theorem 1 implies that $\mathcal{C} = \{\Delta_{A_1, B_1}, \dots, \Delta_{A_s, B_s}\}$ is a maximal collection of pairwise quasi-commuting quantum minors in $\mathbb{C}_q[Mat_{k,m}]$ if and only if $\{S(A_1, B_1), \dots, S(A_s, B_s)\} \sqcup \{[1 \dots k]\}$ is a maximal collection of pairwise weakly separated k -subsets of $[1 \dots k + m]$. This identification is a central component in our attempt to resolve Problem 2. Theorem 1.3 of [7] asserts that the size of any maximal collection of pairwise weakly separated k -subsets of $[1 \dots n]$ is sharply bounded by $k(n - k) + 1$. Setting $n = k + m$ we obtain:

Proposition 1. *The size of any maximal collection of pairwise quasi-commuting quantum minors in $\mathbb{C}_q[Mat_{k \times m}]$ is sharply bounded by km .*

In [7] the following *purity* property is conjectured: all maximal collections of pairwise weakly separated subsets (not necessarily k -subsets) of $[1 \dots n]$ have size $\binom{n+1}{2} + 1$. The analogue of this purity conjecture for k -subsets is given by:

Conjecture 1 (Purity). *All maximal collections of pairwise weakly separated k -subsets of $[1 \dots n]$ have size $k(n - k) + 1$. Equivalently, all maximal collections of pairwise quasi-commuting quantum minors in $\mathbb{C}_q[Mat_{k \times m}]$ have size km .*

In Sections 5 and 6 we prove this assertion for the cases $k = 2$ and $k = 3$ respectively.

In Section 3 we expose a new feature specific to the quantum Grassmannian: quasi-commutativity of the quantum Plücker coordinates in $\mathbb{C}_q[\mathbb{G}_{k,n}]$ is preserved under the natural action of the dihedral group D_n . More precisely, we show that the

natural D_n -action on k -subsets of $[1 \dots n]$ preserves weak separability. We do not know of an analogue of this action for the full quantum flag variety. Let $W(k, n)$ be the set of all maximal collections of pairwise weakly separated k -subsets in $[1 \dots n]$. The induced D_n -action on $W(k, n)$ is instrumental in proving several assertions in this paper.

For a set I and elements x and y let I_{xy} denote $I \cup \{x, y\}$. The set $W(k, n)$ possesses the following interesting structure.

Theorem 3. *Let \mathcal{C} be a maximal collection of pairwise weakly separated k -subsets of $[1 \dots n]$. Suppose that $I_{ij}, I_{it}, I_{js}, I_{st} \in \mathcal{C}$ for some $i < s < j < t$ and for some $I \subset [1 \dots n] - \{i, j, s, t\}$ with $|I| = k - 2$. Then \mathcal{C} contains either I_{ij} or I_{st} and not both. Moreover, the transformation*

$$(1) \quad \mathcal{C} \mapsto \begin{cases} \mathcal{C} - \{I_{ij}\} \sqcup \{I_{st}\} & \text{if } I_{ij} \in \mathcal{C} \\ \mathcal{C} - \{I_{st}\} \sqcup \{I_{ij}\} & \text{if } I_{st} \in \mathcal{C} \end{cases}$$

preserves weak separability and maximality.

This transformation is an analogue of the Yang-Baxter “flip” introduced in [7]; here we refer to these transformations as $(2, 4)$ -moves due to the fact that they originate on $\mathbb{C}_q[\mathbb{G}_{2,4}]$.

Conjecture 2 (Transitivity). *Let \mathcal{C} and \mathcal{B} be any collections in $W(k, n)$. Then there is a sequence of $(2, 4)$ -moves transforming \mathcal{C} into \mathcal{B} .*

If true the conjecture effectively settles Problem 2. In addition it provides a method to obtain all collections in $W(k, n)$: simply propagate a given maximal collection by all possible $(2, 4)$ -moves. In Section 3 we explain why the validity of Conjecture 2 implies the validity of Conjecture 1. In Sections 5 and 6 we prove this Conjecture 2 for the cases $k = 2$ and $k = 3$. In Section 8 we explore applications of this conjecture to total positivity.

In Section 4 we describe certain maximal collections in $W(k, n)$ arising from *double wiring arrangements*. In Section 7 we present a construction that recursively generates all collections in $W(3, n)$ by lifting collections from $W(3, n - 1)$. In principle this construction should provide a method to compute the size of $W(3, n)$.

2. THE QUANTUM GRASSMANNIAN AND PROOFS OF THEOREMS 1 AND 2

Definition 3. *The quantum Grassmannian $\mathbb{C}_q[\mathbb{G}_{k,n}]$, as defined in [10], is the $\mathbb{C}(q)$ -algebra with unity generated by all quantum Plücker coordinates Δ^K where K is a k -subset of $[1 \dots n]$ subject to the relations:*

$$\sum_{i \in I - J} (-q)^{\text{inv}(i, I) - \text{inv}(i, J)} \Delta^{I - \{i\}} \Delta^{J \sqcup \{i\}} = 0$$

for any $(k+1)$ -subset I and $(k-1)$ -subset J . Here $\text{inv}(i, X)$ is the number of $x \in X$ such that $i > x$.

Proposition 2 (Quantum Stieffel-Plücker Correspondence). *There exists a unique $\mathbb{C}(q)$ -algebra embedding $\varphi : \mathbb{C}_q[\text{Mat}_{k \times m}] \longrightarrow \mathbb{C}_q[\mathbb{G}_{k, k+m}]$ such that*

$$\Delta_{I,J} \longmapsto q^{\binom{l}{2}} \Delta^{l-1} \Delta^{S(I,J)}$$

where $l = |I| = |J|$ and $\Delta = \Delta^{[1 \dots k]}$.

Proof. The proof that the Faddeev-Reshetikhin-Takhtadzhyan relations are preserved under the correspondence $x_{i,j} \longmapsto \Delta^{S(\{i\}, \{j\})}$ and that $\Delta_{I,J}$ is sent to $q^{\binom{l}{2}} \Delta^{l-1} \Delta^{S(I,J)}$ is a simple modification of the proof of the quantum analogue of Bazin's theorem presented in Theorem 3.8 of [6].

The classical analogue of φ , obtained by specializing q to 1, is easily seen to be injective. This taken together with Theorem 3.5(c) of [10] and the fact that the monomials consisting of products of lexicographically ordered generators $x_{i,j}$ form a basis for $\mathbb{C}_q[\text{Mat}_{k \times m}]$ over $\mathbb{C}(q)$ proves injectivity of φ . □

It is well known that $\Delta^{[1 \dots k]}$ is quasi-central. Thus Proposition 2 tells us that two quantum minors $\Delta_{A,B}$ and $\Delta_{C,D}$ will quasi-commute exactly when the corresponding quantum Plücker coordinates $\Delta^{S(A,B)}$ and $\Delta^{S(C,D)}$ quasi-commute. In turn, the conditions for two quantum Plücker coordinates to quasi-commute are explained by the following proposition of [7]:

Proposition 3. *Two quantum Plücker coordinates Δ^I and Δ^J in $\mathbb{C}_q[\mathbb{G}_{k,n}]$ quasi-commute if and only if I and J are weakly separated. If I and J satisfy case 1 of Definition 2 then $c(\Delta^I \mid \Delta^J) = |J''| - |J'|$.*

Theorem 1 now follows from Propositions 2 and 3. Theorem 2 also follows from Propositions 2 and 3 along with the fact that $c(\Delta^{|A|-1} \mid \Delta^{S(C,D)}) = |C|(|A| - 1)$ and $c(\Delta^{S(A,B)} \mid \Delta^{|C|-1}) = |A|(1 - |C|)$.

3. PROOF OF THEOREM 3

It is convenient to visualize a k -subset of $[1 \dots n]$ as a subpolygon of the regular polygon with n vertices labeled counter-clockwise by the indices $[1 \dots n]$. Represent the dihedral group D_n as the group of symmetries of the n -gon. Clearly D_n acts on the set of k -subsets of $[1 \dots n]$ under this realization.

Proposition 4. *If two k -subsets I and J of $[1 \dots n]$ are weakly separated then $g(I)$ and $g(J)$ are weakly separated for any $g \in D_n$.*

Proof. In [7] it is shown that I and J are weakly separated precisely when, after interchanging I and J if necessary, either:

- a) $|I| < |J|$ and there do not exist three indices $a < b < c$ such that $I \cap \{a, b, c\} = \{b\}$ and $J \cap \{a, b, c\} = \{a, c\}$ **or**
- b) $|I| = |J|$ and there do not exist four indices $a < b < c < d$ such that $I \cap \{a, b, c, d\} = \{a, c\}$ and $J \cap \{a, b, c, d\} = \{b, d\}$

Part b) above indicates that two k -subsets I and J are weakly separated precisely, when viewed as subpolygons, no diagonal of the subpolygon I disjoint from J crosses a diagonal of J disjoint from I . This property is clearly preserved under any dihedral symmetry of the n -gon.

□

A k -subset I is called *boundary* if it consists of k consecutive indices of the n -gon; i.e. any k -subset of the form $g([1 \dots k])$ for $g \in D_n$. Since $[1 \dots k]$ is weakly separated with every k -subset it follows that the set of all k -boundary subsets is common to every maximal collection of pairwise weakly separated k -subsets.

Proof of Theorem 3:

To prove the first part of the theorem notice that since I_{ij} and I_{st} are not weakly separated it is clear that both can not be in \mathcal{C} . So we need only demonstrate that one of them is present in \mathcal{C} . Given a k -subset J of $[1 \dots n]$ such that J is weakly separated from $I_{is}, I_{sj}, I_{jt}, I_{it}$ and different from I_{ij} and I_{st} we need to show that J is weakly separated from both I_{ij} and I_{st} .

Proposition 4 shows that we may reduce the proof to the case of $t = n$ after suitably translating the collection \mathcal{C} by the *dihedral action*. Assume that $t = n$. Let $J^- = J - \{n\}$. Since $|J| = k$ and J is different from I_{ij} and I_{st} , it follows that J^- is different from both I_{ij} and I_s . By Lemma 3.2 of [7], J^- is weakly separated from I_{is}, I_{sj}, I_j, I_i . By Lemma 5.2 of [7], it follows that J^- is weakly separated from both I_{ij} and I_s and, after an easy application of part b) above, that J is weakly separated from both I_{ij} and I_{sn} , as claimed.

The above argument also shows that the transformation (1) preserves weak separability and maximality, thus concluding the proof of Theorem 3. □

Returning to Conjecture 2, notice that if it is true and if we can find a collection \mathcal{A} in $W(k, n)$ for which $|\mathcal{A}| = k(n - k) + 1$ then Conjecture 1 will follow. One can easily verify that the collection $\mathcal{A} = \mathcal{A}_n$ whose non-boundary sets are

$$\left\{ [1 \dots i] \sqcup [j \dots k + j - i - 1] \mid 1 \leq i < k \text{ and } i + 1 < j < n + i - k \right\}$$

has the desired properties.

4. WIRING ARRANGEMENTS

In [7] a recursive procedure is described through which all maximal families of pairwise weakly separated subsets (not necessarily k -subsets) of $[1 \dots n]$ are obtained. In principle this recursion can be restricted to produce all families in $W(k, n)$. Nevertheless, the process is not very practical. In this section we explore a non-recursive combinatorial device which parametrizes a large portion of the collections in $W(k, n)$. This device is a modification of a construction in [3].

Recall that the symmetric group S_n is generated by the simple reflections $s_i = (i, i+1)$ satisfying the Coxeter relations. A reduced word for an element $g \in S_n$ is sequence of indices i_1, \dots, i_l such that $g = s_{i_1} \cdots s_{i_l}$ with l minimal. For the group $S_k \times S_m$ we will use the indices $[\overline{1}, \dots, \overline{k-1}]$ to label the simple reflections corresponding to the S_k component and the indices $[1 \dots m-1]$ to label the simple reflections for the S_m component. Under this convention a reduced word for an element $(u, v) \in S_k \times S_m$ can be identified with a shuffle of a reduced word for u , written with indices in $[\overline{1}, \dots, \overline{k-1}]$, and a reduced word for v written with indices $[1 \dots m-1]$.

Let $w_0^{(k)}$ and $w_0^{(m)}$ denote the longest elements in S_k and S_m respectively. We say a reduced word for $(w_0^{(k)}, w_0^{(m)}) \in S_k \times S_m$ is *optimal* if the associated reduced word for $w_0^{(m)} \in S_m$ has a total of only $\binom{m-k}{2}$ occurrences of the indices $[k+1, \dots, m-1]$. Given an optimal reduced word \mathbf{i} of $(w_0^{(k)}, w_0^{(m)}) \in S_k \times S_m$ we will manufacture a maximal collection $\mathcal{C}(\mathbf{i})$ of pairwise quasi-commuting quantum minors. This collection is obtained by means of the *double wiring arrangement* $\text{Arr}(\mathbf{i})$ attached to \mathbf{i} , as introduced in [3].

Recall first the definition of a *single wiring arrangement* attached to a reduced word. It is easiest to understand this definition with an example. Consider the reduced word 1231 of the permutation $\begin{pmatrix} 1 & 2 & 3 & 4 \\ 3 & 2 & 4 & 1 \end{pmatrix} \in S_4$. The corresponding single wiring arrangement is:

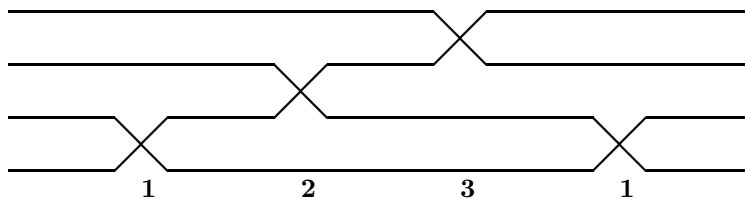


FIGURE 1. Single wiring arrangement

We associate a crossing at the i th level (counting from the bottom up) for each i in the reduced word. To obtain the double wiring arrangement for $(u, v) \in S_k \times S_m$ we superimpose the single wiring arrangements for the reduced words of u and v respectively aligning them closely in the vertical direction (starting at the bottom)

$\mathbf{i} = 2 \, \overline{1} \, 1 \, 2 \, 3 \, \overline{2} \, 2 \, 1 \, 4 \, \overline{1} \, 3 \, 2 \, 1$ for $(w_0^{(3)}, w_0^{(5)}) \in S_3 \times S_5$ is:

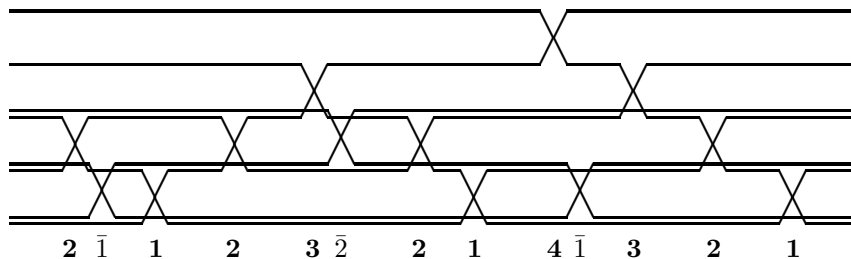


FIGURE 2. Double wiring arrangement

FIGURE 3. Labeled arrangement

$$\text{Let } \mathcal{C}(\mathbf{i}) = \left\{ \Delta_{I(C), J(C)} \mid C \text{ a chamber of } \text{Arr}(\mathbf{i}) \text{ of level } \leq k \right\}.$$

Lemma 1. *Let \mathbf{i} be an optimal reduced word for $(w_0^{(k)}, w_0^{(m)}) \in S_k \times S_m$. Then the size of $\mathcal{C}(\mathbf{i})$ is km .*

Proof. Given \mathbf{i} , the number of chambers in the first k strips of the corresponding double wiring arrangement is equal to the number of red and black crossings in the first k strips plus k - corresponding to the k far right chambers. The number of black (respectively red) crossings in the first k strips in turn is given by the number of simple reflections j (respectively \bar{j}) occurring in the reduced word \mathbf{i} with $1 \leq j \leq k$. The number of \bar{j} in \mathbf{i} with $1 \leq j \leq k$ is $\binom{k}{2}$. The number of j in \mathbf{i} with $1 \leq j \leq k$ is $\binom{m}{2} - \# \{ j \text{ occurring in } \mathbf{i} \mid k+1 \leq j \leq m-1 \}$; if \mathbf{i} is *optimal*

this will be $\binom{m}{2} - \binom{m-k}{2}$. Consequently the number of chambers occurring in the first k strips of the double wiring arrangement for \mathbf{i} optimal (or equivalently the size of $\mathcal{C}(\mathbf{i})$) is:

$$\binom{k}{2} + \binom{m}{2} - \binom{m-k}{2} + k = mk$$

□

Proposition 5. *If \mathbf{i} is an optimal reduced word for $(w_0^{(k)}, w_0^{(m)}) \in S_k \times S_m$ then $\mathcal{C}(\mathbf{i})$ is a maximal collection of pairwise quasi-commuting quantum minors in $\mathbb{C}_q[Mat_{k \times m}]$. Moreover, given $\Delta_{A,B}$ and $\Delta_{I,J}$ in $\mathcal{C}(\mathbf{i})$ either*

$$(2) \quad A - I \prec I - A \quad \text{and} \quad J - B \prec B - J \quad \text{or}$$

$$(3) \quad I - A \prec A - I \quad \text{and} \quad B - J \prec J - B$$

Proof. Take any quantum minors $\Delta_{A,B}$ and $\Delta_{I,J}$ in $\mathcal{C}(\mathbf{i})$. Lemma 4.1 of [7] proves that if X and Y are chamber sets of a single wiring arrangement then either $X - Y \prec Y - X$ or $Y - X \prec X - Y$. This, taken together with the fact that the single wiring arrangements for the S_k and S_m components of \mathbf{i} are oppositely labeled, proves the second part of the proposition.

To prove that $\Delta_{A,B}$ and $\Delta_{I,J}$ quasi-commute we must show that $S(A, B)$ and $S(I, J)$ are weakly separated. We may assume, after exchanging A with I and B with J if necessary, that $A - I \prec I - A$ and $J - B \prec B - J$. This in turn is equivalent to

$$\left(S(A, B) - S(I, J) \right) \cap [1 \dots k] \prec S(I, J) - S(A, B) \prec \left(S(A, B) - S(I, J) \right) - [1 \dots k]$$

which demonstrates that $S(A, B)$ and $S(I, J)$ are weakly separated. The fact that $\mathcal{C}(\mathbf{i})$ is maximal follows from Lemma 1 and Proposition 1.

□

It is possible to prove the converse of Proposition 5, namely: If \mathcal{C} is a collection of quantum minors $\Delta_{A,B}$ whose indices pairwise satisfy either condition 2 or 3, and if \mathcal{C} is maximal with respect to this property, then \mathcal{C} is of the form $\mathcal{C}(\mathbf{i})$ for some optimal reduced word \mathbf{i} .

Given an optimal reduced word \mathbf{i} the following collection is in $W(k, k+m)$:

$$\left\{ S(I(C), J(C)) \parallel C \text{ a chamber of } \text{Arr}(\mathbf{i}) \text{ of level } \leq k \right\} \sqcup \left\{ [1 \dots k] \right\}$$

In the case of $W(3, 6)$ all collections are obtained via double wiring arrangements. There are 34 in total and they are explicitly described in [3] and [4]. Every maximal family in $W(3, 6)$ is dihedrally equivalent to one of the following five collections (we omit boundary sets):

$$\begin{aligned} & \left\{ \{124\}, \{125\}, \{134\}, \{145\} \right\} \quad \left\{ \{124\}, \{125\}, \{145\}, \{245\} \right\} \\ & \left\{ \{124\}, \{134\}, \{145\}, \{146\} \right\} \quad \left\{ \{125\}, \{134\}, \{135\}, \{145\} \right\} \\ & \left\{ \{135\}, \{136\}, \{145\}, \{235\} \right\} \end{aligned}$$

In general it is not the case that every maximal collection in $W(k, n)$ corresponds to some double wiring arrangement, even after dihedral translation. This is evidenced already in the case of $\mathbb{C}_q[\mathbb{G}_{2,n}]$. In Section 5 we shall demonstrate such a maximal collection.

5. THE CASE OF $\mathbb{C}_q[\mathbb{G}_{2,n}]$

We identify the 2-subsets of $[1 \dots n]$ with chords inscribed in a regular n -gon. Clearly two 2-subsets of $[1 \dots n]$ are weakly separated if and only if the corresponding chords do not cross in the interior of the polygon. Under this identification collections $\mathcal{C} \in W(2, n)$ correspond to maximal collections of non-crossing chords - i.e. triangulations of an n -gon.

Theorem 4 (Transitivity). *Let $\mathcal{C}, \mathcal{B} \in W(2, n)$. Then there is a sequence of $(2, 4)$ -moves transforming \mathcal{C} into \mathcal{B} .*

Proof. This theorem follows from the well known fact that the any two triangulations are connected by a series of chord exchanges where the diagonal chord of an inscribed quadrilateral is "flipped" to its crossing pair. The diagonal "flips" correspond to $(2, 4)$ -moves. \square

Corollary 1 (Purity). *Let $\mathcal{C} \in W(2, n)$. Then $|\mathcal{C}| = 2(n - 2) + 1$.*

Proof. Immediate corollary of Theorem 4. \square

Since $W(2, n)$ is identified with the set of triangulations of an n -gon it follows that $|W(2, n)|$ is the Catalan number $\frac{1}{n-1} \binom{2n-4}{n-2}$. For $k > 2$ the size of $W(k, n)$ is not known.

In [5] it is shown that the coordinate ring $\mathbb{C}[\mathbb{G}_{2 \times n}]$ has a basis consisting of all monomials of Plücker coordinates whose indices are pairwise weakly separated. Using the quantum short Plücker relation given by

$$\Delta^{Iij} \Delta^{Ist} = q \Delta^{Iis} \Delta^{Ijt} + q^{-1} \Delta^{Iit} \Delta^{Isj}$$

for $i < s < j < t$ as a straightening rule, we obtain the following quantum analogue of this result:

Proposition 6. *The set of all monomials consisting of lexicographically ordered pairwise quasi-commuting quantum Plücker coordinates is a basis for $\mathbb{C}_q[\mathbb{G}_{2,n}]$.*

Using Proposition 5 and the identification of maximal collections in $W(2, n)$ with triangulations of an n -gon we can characterize those maximal collections which can be parametrized, up to the dihedral action, by double wiring arrangements. Given $\mathcal{C} \in W(2, n)$ there exists $g \in D_n$ for which $g \cdot \mathcal{C}$ is parametrized by a double wiring arrangement if and only if there exists an external edge of the polygon (i.e. a boundary 2-set) such that for any other external edge there is **no** chord in the associated triangulation, which separates both the edges and is disjoint from both. The following collection in $W(2, 9)$, represented as a triangulation, is an example of a collection which is not parametrized, up to the dihedral action, by a double wiring arrangement:

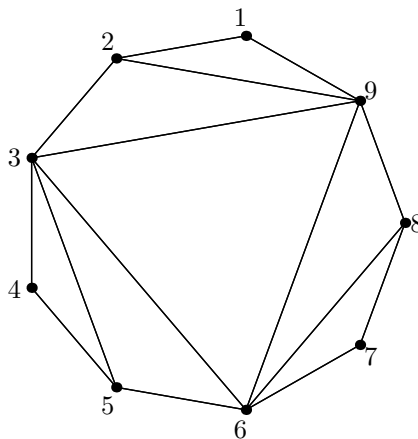


FIGURE 4. Non-Parametrized $W(2, 9)$ collection

6. THE CASE OF $\mathbb{C}_q[\mathbb{G}_{3,n}]$

In this section we prove the Transitivity and Purity Conjectures for $k = 3$.

Theorem 5 (Transitivity). *Let $\mathcal{C}, \mathcal{B} \in W(3, n)$. Then there is a sequence of $(2, 4)$ -moves transforming \mathcal{C} into \mathcal{B} .*

Corollary 2 (Purity). *Let $\mathcal{C} \in W(3, n)$ then $|\mathcal{C}| = 3(n - 3) + 1$.*

Proof of Transitivity:

The essential strategy is to show that any collection $\mathcal{C} \in W(3, n)$ can be reduced by a sequence of $(2, 4)$ -moves to the "base" collection \mathcal{A}_n whose non-boundary 3-sets are

$$\left\{ \{1, s, s+1\} \parallel 2 < s < n-1 \right\} \sqcup \left\{ \{1, 2, s\} \parallel 3 < s < n \right\}$$

We first prove that whenever a collection \mathcal{C} can be $(2, 4)$ -reduced to \mathcal{A}_n then so can any of its dihedral translations $g \cdot \mathcal{C}$ for $g \in D_n$. In Lemma 3 we then show that any maximal collection can be translated dihedrally to a maximal collection containing the 3-set $\{1, n-2, n-1\}$. We conclude the proof by showing that any such collection can be reduced by a sequence of $(2, 4)$ -moves to the collection \mathcal{A}_n .

Lemma 2. *Let $\mathcal{C} \in W(3, n)$. If \mathcal{C} can be reduced by a sequence of $(2, 4)$ -moves to \mathcal{A}_n then so can the collection $g \cdot \mathcal{C}$ for any $g \in D_n$.*

Proof. Since the D_n -action preserves $(2, 4)$ -moves it is enough to verify this assertion in the case where $\mathcal{C} = \mathcal{A}_n$.

Proceed by induction on n . For $n \leq 4$ the statement is evident. Assume $n > 4$. It is enough to verify the claim for the group elements ρ_n and σ_n , which generate D_n , given by

$$\rho_n = \begin{pmatrix} 1 & 2 & \cdots & n-1 & n \\ 2 & 3 & \cdots & n & 1 \end{pmatrix} \quad \sigma_n = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & \cdots \\ 2 & 1 & n & n-1 & n-2 & \cdots \end{pmatrix}$$

This follows from the observation that if $g \cdot \mathcal{C}$ can be reduced by a sequence of $(2, 4)$ -moves to \mathcal{B} then $hg \cdot \mathcal{C}$ can be reduced to $h \cdot \mathcal{B}$.

The collection $\sigma_n \cdot \mathcal{A}_n$ contains the 3-sets $\{1, 2, n-1\}$, $\{2, n-2, n-1\}$, $\{n-2, n-1, n\}$, $\{1, n-1, n\}$, and $\{2, n-1, n\}$. Applying the $(2, 4)$ -move which replaces $\{2, n-1, n\}$ with $\{1, n-2, n-1\}$ we obtain $\sigma_{n-1} \cdot \mathcal{A}_{n-1} \sqcup \left\{ \{1, 2, n\}, \{1, n-1, n\}, \{n-2, n-1, n\} \right\}$. By induction $\sigma_{n-1} \cdot \mathcal{A}_{n-1}$ can be reduced by a sequence of $(2, 4)$ -moves to \mathcal{A}_{n-1} . Thus $\sigma_n \cdot \mathcal{A}_n$ can be reduced to $\mathcal{A}_{n-1} \sqcup \left\{ \{1, 2, n\}, \{1, n-1, n\}, \{n-2, n-1, n\} \right\} = \mathcal{A}_n$.

To deal with ρ_n , notice that $\rho_n \cdot \mathcal{A}_n$ contains the 3-sets $\{1, 2, n\}$, $\{1, 2, 3\}$, $\{2, 3, n-1\}$, $\{2, n-1, n\}$, and $\{2, 3, n\}$. We apply the $(2, 4)$ -move which replaces $\{2, 3, n\}$ with $\{1, 2, n-1\}$. This new collection contains the 3-sets $\{1, n-1, n\}$, $\{1, 2, n-1\}$, $\{2, n-2, n-1\}$, $\{n-2, n-1, n\}$, and $\{2, n-1, n\}$. We may apply the $(2, 4)$ -move which replaces $\{2, n-1, n\}$ with $\{1, n-1, n-2\}$. The resulting collection is exactly $\rho_{n-1} \cdot \mathcal{A}_{n-1} \sqcup \left\{ \{1, 2, n\}, \{1, n-1, n\}, \{n-2, n-1, n\} \right\}$. By the induction hypothesis $\rho_{n-1} \cdot \mathcal{A}_{n-1}$ can be reduced by a sequence of $(2, 4)$ -moves to \mathcal{A}_{n-1} . Consequently $\rho_n \cdot \mathcal{A}_n$ can be reduced to $\mathcal{A}_{n-1} \sqcup \left\{ \{1, 2, n\}, \{1, n-1, n\}, \{n-2, n-1, n\} \right\} = \mathcal{A}_n$.

□

Lemma 3. *Given $\mathcal{C} \in W(3, n)$ there exists $g \in D_n$ such that $g \cdot \mathcal{C}$ contains the 3-set $\{1, n-2, n-1\}$.*

Proof. For a 3-subset I of $[1 \dots n]$ define the diameter of I to be the minimal cardinality of a boundary k -subset of $[1 \dots n]$ that contains I . Thus the boundary 3-subsets are precisely those of diameter 3. Let us call 3-subsets of diameter 4 *almost boundary* subsets. It suffices to prove that every maximal collection \mathcal{C} contains an almost boundary subset.

Assume by contradiction that \mathcal{C} does not contain an almost boundary 3-subset. We make the following easy observation:

Remark 1. *Let a, b, c , and d be four consecutive vertices in $[1 \dots n]$; then the 3-subsets that are not weakly separated with an almost boundary subset $\{a, c, d\}$ are precisely the non-boundary 3-subsets containing b but not a .*

Therefore our assumption and maximality of \mathcal{C} imply that for every two consecutive vertices a and b in $[1 \dots n]$, there is a non-boundary 3-subset in \mathcal{C} which contains b but not a .

Choose a non-boundary 3-subset $\{a, c, d\}$ in \mathcal{C} of minimal possible diameter. Without loss of generality, we can assume that a boundary subset of minimal cardinality that contains $\{a, c, d\}$ has a and d as its endpoints; let us denote this boundary subset by $[a, d]$. We can also assume that c is not a neighbor of a . Let b be the neighbor of a in $[a, d]$. Consider a 3-subset I in \mathcal{C} such that I contains b but not a . Since I is weakly separated from $\{a, c, d\}$ it must be contained in $[b, d] = [a, d] - \{a\}$. But then I has smaller diameter than $\{a, c, d\}$ which contradicts our choice of $\{a, c, d\}$. This proves the claim and hence the lemma as well. □

For any collection $\mathcal{C} \in W(3, n)$ we define its *height* $H(\mathcal{C})$ to be the number of non-boundary 3-sets containing n . An immediate consequence of Remark 1 is that $H(\mathcal{C}) = 0$ if and only if both $\{1, 2, n-1\}$ and $\{1, n-2, n-1\}$ are in \mathcal{C} .

Lemma 4. *Let $\mathcal{C} \in W(3, n)$ with $\{1, n-2, n-1\} \in \mathcal{C}$. Then \mathcal{C} can be reduced by a sequence of $(2, 4)$ -moves to a collection of height $H = 0$.*

Proof. We proceed by induction on the height. If $H(\mathcal{C}) = 0$ then we are already done. Assume inductively that the assertion is true for collections of height $H = k \geq 0$ and let \mathcal{C} be a collection of height $H(\mathcal{C}) = k + 1$. We need the following:

Lemma 5. *Let $\mathcal{C} \in W(3, n)$ and suppose that $\{1, n-2, n-1\} \in \mathcal{C}$. Then there exists a unique index $b > 1$ such that both $\{1, b, n-1\}$ and $\{1, b, n\}$ are in \mathcal{C} . We call b the *pinch point* over n and $n-1$.*

Proof. Let b be the maximal index with the property that $\{1, b, n\} \in \mathcal{C}$. Suppose, by contradiction, that $\{1, b, n-1\} \notin \mathcal{C}$. By maximality of \mathcal{C} this means there exists a non-boundary set $I \in \mathcal{C}$ which is not weakly separated with $\{1, b, n-1\}$. Therefore there exist indices $s, t \in I$ such that one of the following holds:

1. $1 < s < b < t < n-1$
2. $1 < s < b$ and $t = n$
3. $b < s < n-1$ and $t = n$

Case 1: Since I and $\{1, b, n\}$ are weakly separated it follows that $b \in I$. But then I will be weakly separated with $\{1, b, n-1\}$.

Case 2: Since $\{1, n-2, n-1\} \in \mathcal{C}$ and since I is a non-boundary set containing n it follows that $1 \in I$. But then I will be weakly separated with $\{1, b, n-1\}$.

Case 3: Once again it must be the case that $1 \in I$. So $I = \{1, s, n\}$ where $b < s$ violating the maximality of b .

Hence $\{1, b, n-1\} \in \mathcal{C}$. Suppose there was another pinch point $b' \neq b$. Either $b' < b$ or $b' > b$. If $b' < b$ then $\{1, b', n-1\}$ will not be weakly separated from $\{1, b, n\}$. If $b' > b$ then $\{1, b', n\}$ will not be weakly separated from $\{1, b, n-1\}$. Both possibilities violate that fact that \mathcal{C} consists of only pairwise weakly separated 3-sets. Uniqueness follows. \square

Lemma 6. *Let $\mathcal{C} \in \mathcal{W}(3, n)$ and assume $\{1, n-2, n-1\} \in \mathcal{C}$. Let b be the pinch point over n and $n-1$. Assume in addition that $b > 2$. Then there exists a with $1 < a < b$ such that both $\{1, a, b\}$ and $\{1, a, n\}$ are in \mathcal{C} .*

Proof. Consider the set of all x with the property that $x < b$ and $\{1, x, n\} \in \mathcal{C}$. This set is clearly non-empty since $2 < b$ and $\{1, 2, n\} \in \mathcal{C}$. Let a be the maximal index with this property. Suppose $\{1, a, b\} \notin \mathcal{C}$. Then there exists $I \in \mathcal{C}$ with $s, t \in I$ such that one of the following holds:

1. $1 < s < a < t < b$
2. $1 < s < a < b < t$
3. $a < s < b < t$

Case 1: Since $\{1, a, n\} \in \mathcal{C}$ it follows that I and $\{1, a, n\}$ must be weakly separated. The only way this can happen is that $a \in I$. But then I and $\{1, a, b\}$ will be weakly separated.

Case 2: Since I and $\{1, a, n\}$ are weakly separated it must be the case that $t = n$. Since $\{1, n-2, n-1\} \in \mathcal{C}$ it follows that I and $\{1, n-2, n-1\}$ are weakly separated. The only way this can be resolved is that $1 \in I$. But then I and $\{1, a, b\}$ are weakly separated.

Case 3: Either $t = n$ or not. Suppose $t \neq n$. Since $\{1, b, n\} \in \mathcal{C}$, and hence weakly separated from I , it follows that $b \in I$ in which case I and $\{1, a, b\}$ will be weakly separated. Thus $t = n$. Since $\{1, b, n-1\} \in \mathcal{C}$ we know that I and $\{1, b, n-1\}$ are weakly separated. The only way this can happen is that $1 \in I$ and hence $I = \{1, s, n\}$. But this violates the maximality of a since $a < s < b$.

Thus $\{1, a, b\}$ and $\{1, a, n\}$ are in \mathcal{C} as required. \square

Returning to Lemma 4, let b be the pinch point of \mathcal{C} - i.e. the unique index b such that both $\{1, b, n-1\}$ and $\{1, b, n\}$ are in \mathcal{C} . If $b = 2$ it follows that $\{1, 2, n-1\} \in \mathcal{C}$.

This, taken together with the fact that $\{1, n-2, n-1\} \in \mathcal{C}$, violates the hypothesis that $H(\mathcal{C}) > 0$. Therefore $b > 2$.

Since $b > 2$ Lemma 6 implies that there exists a with $1 < a < b$ such that both $\{1, a, b\}$ and $\{1, a, n\}$ are in \mathcal{C} . Thus \mathcal{C} contains $\{1, a, b\}$, $\{1, a, n\}$, $\{1, b, n-1\}$, $\{1, b, n\}$, and $\{1, n-1, n\}$. The associated $(2, 4)$ -move for this quintuple replaces $\{1, b, n\}$ with $\{1, a, n-1\}$. Let \mathcal{B} be the resulting collection. Notice that \mathcal{B} contains $\{1, n-2, n-1\}$ and that $H(\mathcal{B}) = H(\mathcal{C}) - 1 = k$. By induction \mathcal{B} can be further reduced by a sequence of $(2, 4)$ -moves into a collection of height $H = 0$. Concantening this $(2, 4)$ -reduction with the $(2, 4)$ -move transforming \mathcal{C} to \mathcal{B} we obtain the desired reduction for \mathcal{C} .

□

Now we are ready to finish the proof of Transitivity. Let $\mathcal{C} \in W(3, n)$. By Lemma 3 there is $g \in D_n$ such that $g \cdot \mathcal{C}$ contains the 3-set $\{1, n-2, n-1\}$. By Lemma 4 the collection $g \cdot \mathcal{C}$ can be reduced by a sequence of $(2, 4)$ -moves to a collection \mathcal{B} with height $H(\mathcal{B}) = 0$. The collection $\mathcal{B} - \left\{ \{1, 2, n\}, \{1, n-1, n\}, \{n-2, n-1, n\} \right\}$ is in $W(3, n-1)$ and by induction on n we can assume that it can be reduced by a sequence of $(2, 4)$ -moves to \mathcal{A}_{n-1} . Equivalently \mathcal{B} can be reduced by a sequence of $(2, 4)$ -moves to \mathcal{A}_n . Consequently $g \cdot \mathcal{C}$ can be reduced to \mathcal{A}_n and applying Lemma 2 we conclude that \mathcal{C} can be reduced to \mathcal{A}_n as required.

7. REDUCTION

In this section we present a recursive procedure to generate collections in $W(3, n)$.

Given a 3-subset I of $[1 \dots n]$, we define

$$I' = \begin{cases} I \sqcup \{n-1\} - \{n\} & \text{if } n \in I \text{ and } n-1 \notin I \\ \phi & \text{if } n \in I \text{ and } n-1 \in I \\ I & \text{if } n \notin I \end{cases}$$

For $\mathcal{C} \in W(3, n)$ let $\mathcal{C}' = \{ I' \mid I \in \mathcal{C} \}$, and define $F_{\mathcal{C}}$ to be the set of indices $b \in [2 \dots n-1]$ with $\{1, b, n\} \in \mathcal{C}$ such that $\{1, b\} - \{s, t\} \prec \{s, t\} - \{1, b\}$ whenever $\{s, t, n\} \in \mathcal{C}$ for $1 < s < t$. If \mathcal{C} contains $\{1, n-2, n-1\}$, let $b_{\mathcal{C}}$ be the pinch point of \mathcal{C} (see Lemma 5), that is, the unique index such that both $\{1, b_{\mathcal{C}}, n-1\}$ and $\{1, b_{\mathcal{C}}, n\}$ are in \mathcal{C} .

Theorem 6 (Reduction). *Let $n \geq 4$. The mapping $\mathcal{C} \mapsto (\mathcal{C}', b_{\mathcal{C}})$ defines a bijection between collections in $W(3, n)$ containing $\{1, n-2, n-1\}$ and the set*

$$\left\{ (\mathcal{B}, b) \in W(3, n-1) \times [2 \dots n-2] \mid b \in F_{\mathcal{B}} \right\}$$

The inverse bijection sends a pair (\mathcal{B}, b) to the collection $\hat{\mathcal{B}}_b := \{ I_b \parallel I \in \mathcal{B} \} \sqcup \{ \{1, b, n-1\}, \{1, n-1, n\}, \{n-2, n-1, n\} \}$ where

$$I_b = \begin{cases} I - \{n-1\} \sqcup \{n\} & \text{if } n-1 \in I \text{ and } I - \{1, b, n-1\} \prec \{1, b\} - I \\ I & \text{otherwise} \end{cases}$$

Since by Lemma 3, every collection in $W(3, n)$ is dihedrally equivalent to one containing the near boundary subset $\{1, n-2, n-1\}$, it follows from Theorem 6 that all collections in $W(3, n)$ can be obtained by first lifting collections in $W(3, n-1)$ by the inverse of the reduction procedure and then translating them suitably by the dihedral action.

Proof of Reduction Theorem:

The following lemma shows that the mapping $\mathcal{C} \mapsto (\mathcal{C}', b_c)$ is well defined.

Lemma 7. *Let $\mathcal{C} \in W(3, n)$. Then $\mathcal{C}' \in W(3, n-1)$, and $b_c \in F_{\mathcal{C}'}$.*

Proof. Momentary consideration reveals that \mathcal{C}' consists of pairwise weakly separated 3-subsets of $[1 \dots n-1]$. In virtue of Corollary 2 we know that \mathcal{C}' will be maximal if and only if $|\mathcal{C}'| = 3(n-4) + 1$. Since $\{1, n-2, n-1\} \in \mathcal{C}$ it follows that if $I \in \mathcal{C}$ and $I' = \emptyset$ then either $I = \{1, n-1, n\}$ or $I = \{n-2, n-1, n\}$. Consequently $|\mathcal{C}'| \leq |\mathcal{C}| - 2$. For $I, J \in \mathcal{C}$ if $I' = J'$ then either $I = J$ or else there exists $b \in [2 \dots n-2]$ such that, after interchanging I and J if necessary, $I = \{1, b, n-1\}$ and $J = \{1, b, n\}$. By Lemma 5, b is unique. Hence $|\mathcal{C}'| = |\mathcal{C}| - 3 = 3(n-4) + 1$ as required. The inclusion $b_c \in F_{\mathcal{C}'}$ is also clear from the definitions. \square

To prove that the inverse correspondence is well defined, we need to show that $\hat{\mathcal{B}}_b \in W(3, n)$ and $\{1, n-2, n-1\} \in \hat{\mathcal{B}}_b$ for any $\mathcal{B} \in W(3, n-1)$ and $b \in F_{\mathcal{B}}$. Simple consideration shows that all 3-subsets in $\hat{\mathcal{B}}_b$ are weakly separated because $b \in F_{\mathcal{B}}$. Since \mathcal{B} is maximal we know by Corollary 2 that $|\mathcal{B}| = 3(n-4) + 1$ and thus $|\hat{\mathcal{B}}_b| = |\mathcal{B}| + 3 = 3(n-3) + 1$. Corollary 2 implies that $\hat{\mathcal{B}}_b \in W(3, n)$. Notice also that $\{1, n-2, n-1\} \in \hat{\mathcal{B}}_b$ since $b \leq n-2$.

It remains to show that the mappings $\mathcal{C} \mapsto (\mathcal{C}', b_c)$ and $(\mathcal{B}, b) \mapsto \hat{\mathcal{B}}_b$ are inverse to each other. First suppose that $\mathcal{C} = \hat{\mathcal{B}}_b$. Since both $\{1, b, n\}$ and $\{1, b, n-1\}$ are in $\hat{\mathcal{B}}_b$, the desired equality $(\mathcal{C}', b_c) = (\mathcal{B}, b)$ follows from Lemma 5. Finally, the equality $\hat{\mathcal{C}}'_b = \mathcal{C}$ for $b = b_c$ is clear from the definitions. \square

Example: Let \mathcal{C} be the collection in $W(3, 6)$ whose non-boundary 3-sets are

$$\{ \{136\}, \{146\}, \{236\}, \{346\} \}$$

Here $F_c = \{2, 3\}$. Notice that $4 \notin F_c$ because $\{1, 4\} - \{23\} \not\prec \{2, 3\} - \{1, 4\}$. The index 5 is not present for the same reason. The two possible lifts of \mathcal{C} (omitting boundaries) are:

$$\begin{aligned}\hat{\mathcal{C}}_2 &= \left\{ \{126\}, \{136\}, \{146\}, \{156\}, \{236\}, \{346\} \right\} \\ \hat{\mathcal{C}}_3 &= \left\{ \{137\}, \{136\}, \{146\}, \{156\}, \{236\}, \{346\} \right\}\end{aligned}$$

8. POSITIVITY

Let $\mathbb{G}_{k,n}(\mathbb{C})$ be the Grassmannian of k -subspaces in \mathbb{C}^n . Recall that any k -subspace in $\mathbb{G}_{k,n}(\mathbb{C})$ can be represented by a $k \times n$ matrix whose rows span the k -subspace. The Plücker coordinates are the maximal minors of this $k \times n$ matrix. We say a point $p \in \mathbb{G}_{k,n}(\mathbb{C})$ is positive if it can be represented by a $k \times n$ matrix whose Plücker coordinates $\Delta^I(p)$ are positive real numbers.

Definition 4. Let \mathcal{C} be a collection of k -subsets of $[1 \dots n]$. We say that \mathcal{C} is a positivity test if $p \in \mathbb{G}_{k,n}(\mathbb{C})$ is positive if and only if all $\Delta^I(p)$ are real and positive for each $I \in \mathcal{C}$.

In [7] it is conjectured that maximal families of pairwise weakly separated subsets (not necessarily k -subsets) of $[1 \dots n]$ give rise to positivity tests for the flag variety of type A_n . The analogue of this result for the Grassmannian $\mathbb{G}_{k,n}(\mathbb{C})$ is:

Theorem 7. Let $k = 2$ or $k = 3$. If \mathcal{C} is a maximal collection of pairwise weakly separated k -subsets of $[1 \dots n]$ then the associated collection of Plücker coordinates $\{ \Delta^I \mid I \in \mathcal{C} \}$ is a positivity test.

Proof. Let $\mathcal{C} \in W(k, n)$ and suppose that all $\Delta^I(p)$ are real and positive for $I \in \mathcal{C}$. We need to show that all other Plücker coordinates $\Delta^J(p)$ are real and positive. Take any $J \notin \mathcal{C}$. Take any maximal collection \mathcal{B} containing J . Since k is either 2 or 3 we know that Conjecture 2 holds and thus \mathcal{C} and \mathcal{B} are connected by a sequence of $(2, 4)$ -moves.

Claim: Suppose \mathcal{A} is in $W(k, n)$ and is a positivity test. Let \mathcal{B} be in $W(k, n)$ and assume that \mathcal{B} is obtained from \mathcal{A} by a single $(2, 4)$ -move. Then \mathcal{B} is a positivity test.

Indeed, since \mathcal{A} and \mathcal{B} differ by a single $(2, 4)$ -move there exist $i < s < j < t$ and I , where I is empty if $k = 2$ and $|I| = 1$ if $k = 3$, such that I_{is} , I_{sj} , I_{jt} , and I_{it} are in both \mathcal{A} and \mathcal{B} and such that, without loss of generality, \mathcal{B} is obtained from \mathcal{A} by replacing I_{ij} with I_{st} . The fact that \mathcal{B} is a positivity test is an immediate consequence of the short Plücker relation

$$\Delta^{I_{ij}} \Delta^{I_{st}} = \Delta^{I_{is}} \Delta^{I_{jt}} + \Delta^{I_{it}} \Delta^{I_{sj}}$$

Let l be the minimal number of $(2, 4)$ -moves required to join \mathcal{B} and \mathcal{C} . To prove the theorem proceed by induction on l and use the claim.

□

A positivity test \mathcal{C} is **minimal** if it has no proper subset which is also a positivity test. We conjecture that \mathcal{C} is a minimal positivity test for $\mathbb{G}_{k,n}(\mathbb{C})$ if and only if \mathcal{C} is in $W(k, n)$. In addition, A. Zelevinsky and S. Fomin conjecture that collections \mathcal{C} in $W(k, n)$ have the property that any Plücker coordinate Δ^J can be uniquely expressed as a positive Laurent polynomial in the Plücker coordinates Δ^I for $I \in \mathcal{C}$. The author intends to investigate these issues related to positivity in a forthcoming article.

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